

TERMINAL CHORDS IN CONNECTED CHORD DIAGRAMS

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ABSTRACT. Rooted connected chord diagrams form a nice class of combinatorial objects. Recently they were shown to index solutions to certain Dyson-Schwinger equations in quantum field theory. Key to this indexing role are certain special chords which are called terminal chords. Terminal chords provide a number of combinatorially interesting parameters on rooted connected chord diagrams which have not been studied previously. Understanding these parameters better has implications for quantum field theory.

Specifically, we show that the distributions of the number of terminal chords and the number of adjacent terminal chords are asymptotically Gaussian with logarithmic means, and we prove that the average index of the first terminal chord is $2n/3$. Furthermore, we obtain a method to determine any next-toⁱ leading log expansion of the solution to these Dyson-Schwinger equations, and have asymptotic information about the coefficients of the log expansions.

1. INTRODUCTION

In this paper we are interested in looking at the asymptotic behaviour of some rich and interesting, but somewhat unusual parameters on the combinatorial class of rooted connected chord diagrams. Specifically, we are interested in certain chords known as *terminal chords* which form the base case for a recursive decomposition of rooted connected chord diagrams and the indices of the terminal chords in a recursive ordering of the chords. The reason for investigating these parameters is that they arose in [12] in series solutions to certain Dyson-Schwinger equations in quantum field theory. In order to derive meaningful physics from these series solutions we need to better understand the asymptotics of these parameters. The present paper is a first step towards this understanding. Furthermore the combinatorics of these objects is interesting in its own right and these particular parameters are largely uninvestigated so far.

1.1. COMBINATORIAL SETTING

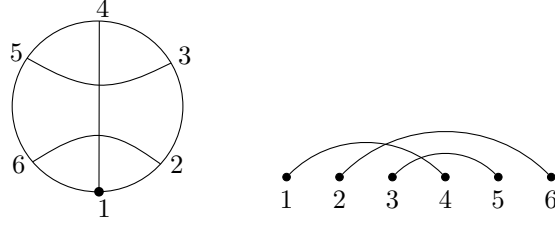
Before explaining the physics context, let us set up what we need for chord diagrams.

Definition 1. A perfect matching of a finite set S is a set of pairs of S such that every element of S is in exactly one pair. A chord diagram with n chords is a perfect matching of $\{1, 2, \dots, 2n\}$. The root chord of a chord diagram is the pair including 1.

As implied by the name, it is convenient to represent chord diagrams with dots and chords. Two conventions coexist in the literature: the circular one and the linear one. They respectively consist in drawing points $1, 2, \dots, 2n$ on a circle in counterclockwise order, or on a line from left to right, and joining by a chord every two elements belonging to the same pair. The matching $\{\{1, 4\}, \{2, 6\}, \{3, 5\}\}$ has been drawn in these two ways in Figure 1. The circular convention has been used in the previous papers, like [12], but we are going to adopt here the linear convention for the rest of the document.

Definition 2. The oriented intersection graph of a chord diagram C is the digraph with a vertex for each chord of C and an oriented edge from chord $\{a, b\}$ to chord $\{c, d\}$ whenever $a < c < b < d$. A

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FIGURE 1. *Left: circular convention. Right: linear convention.*

chord diagram is connected if its oriented intersection graph is connected. A chord is terminal if its vertex in the oriented intersection graph has no outgoing edges.

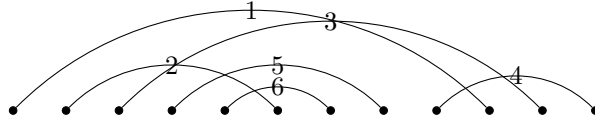
For instance, the oriented intersection graph of the chord diagram of Figure 1 is the tree where $\{1, 4\}$ is the root vertex, and $\{2, 6\}$, $\{3, 5\}$ its two children. This chord diagram is connected, and the terminal chords are $\{2, 6\}$ and $\{3, 5\}$.

The chords inherit an order by the smaller of their endpoints. This is not the order that we want to be working with.

Definition 3. *The intersection order of the chords of a rooted connected chord diagram C is defined as follows.*

- *The root chord of C is the first chord in the intersection order.*
- *Remove the root chord of C and let C_1, C_2, \dots, C_n be the connected components of the result ordered by their first vertex.*
- *For the intersection order of C , after the root chord come all the chords of C_1 ordered inductively in the intersection order, then all the chords of C_2 ordered by intersection order, and so on.*

The chord diagram of Figure 2 is an example of a chord diagram where the intersection order is different from the order by the smaller of their endpoints.

FIGURE 2. *Example of a connected chord diagram and its intersection order.*

Our primary interest is in the terminal chords and their indices in intersection order. We are interested in questions such as

- How many terminal chords does a chord diagram have?
- What is the index of the first terminal chord?
- How many pairs of terminal chords are adjacent in the intersection order?
- What can we say about the gaps between indices of successive terminal chords in intersection order?

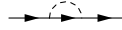
Now we are going to explain why these questions are relevant from a physical point of view.

1.2. PHYSICAL BACKGROUND

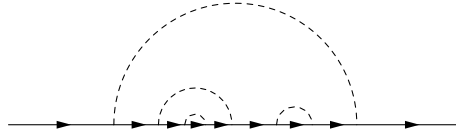
Dyson-Schwinger equations are an important class of equations in quantum field theory. They are the quantum analogues of the classical equations of motion. They are usually written as integral equations and their recursive structure mirrors the decomposition of Feynman graphs into subgraphs.

Because of this recursive structure there is another, more combinatorial way to think about them. Namely, they are functional equations for a sort of weighted generating function. More specifically, the Green functions of a quantum field theory can be thought of as the sum over all Feynman graphs of the theory satisfying certain properties (for example graphs which are 1 particle irreducible – that is 2-edge-connected – and have a fixed set of external edges) and weighted by their Feynman integrals. Thus the Green functions are weighted generating functions of Feynman graphs with highly nontrivial weights. The Green functions are solutions to the Dyson-Schwinger equations, or, viewed the other way around, the Dyson-Schwinger equations are certain functional equations for these weighted generating functions.

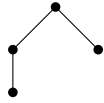
In [12] one of the authors along with Nicolas Marie looked at one particular family of Dyson-Schwinger equations given below in (1). This family of Dyson-Schwinger equations corresponds to the physical situation where we consider all graphs made by inserting a fixed one loop propagator graph into itself in one insertion place. Combinatorially this means that the graphs we are interested in are in bijection with plane rooted trees. For instance, inserting the graph



into itself in all possible ways gives a class of graphs which fits into this situation. One example from this class is



which corresponds to the rooted tree



The Dyson-Schwinger equations considered in [12] are those which can be written in the following form

$$G(x, L) = 1 - xG\left(x, \frac{\partial}{\partial(-\rho)}\right)^{-1} (e^{-L\rho} - 1)F(\rho)|_{\rho=0} \quad (1)$$

where $F(\rho)$ is the Laurent expansion of a regularized Feynman integral for the one loop graph which generated the graph class in question. Note that $G(x, \partial/\partial(-\rho))^{-1}$ is acting as a differential operator on $(e^{-L\rho} - 1)F(\rho)$. Analytically there are subtleties since such an operator is only a pseudo-differential operator. However, we are concerned here solely with series and so interpreting (1) as an equation in formal series everything is well-defined.

For the specific example graphs given above, viewed in Yukawa theory, the equation was solved by Broadhurst and Kreimer in [5]. In this case $F(\rho) = 1/(\rho(1 - \rho))$. They didn't write the Dyson-Schwinger equation in the form of (1) but rather in a more usual physical form as an integral equation (see Examples 3.5 and 3.7 of [18] for how to convert to the form above). However, for the present purposes, we can simply take (1) as the starting point.

The main result of [12] is a series solution to (1) indexed by rooted connected chord diagrams.

Theorem 4 (Theorem 4.13 of [12]). *Suppose $F(\rho) = \sum_{i \geq 0} f_i \rho^{i-1}$. Given a rooted connected chord diagram C let the indices of its terminal chords in intersection order be $b(C) = t_1 < t_2 < \dots < t_k$. Then*

$$G(x, L) = 1 - \sum_{i \geq 1} \frac{(-L)^i}{i!} \sum_{\substack{C \\ b(C) \geq i}} x^{|C|} f_{b(C)-i} f_0^{|C|-k} \prod_{j=2}^k f_{t_j-t_{j-1}} \quad (2)$$

solves (1), where the sum is over rooted connected chord diagrams with the indicated restriction.

Note that k and the t_j depend on C in (2), but this has been left implicit to keep the notation from getting too heavy. More general Dyson-Schwinger equations have similar chord diagram expansions (see [9]).

To better understand $G(x, L)$ we see from (2) that the keys are to understand the number of terminal chords, the index of the first terminal chord, and the differences between indices of successive terminal chords.

1.3. STRUCTURE OF THE DOCUMENT

In Section 2, we set up the enumerative background of the present article. We introduce what is known on exact and asymptotic enumeration of connected chord diagrams. We also give some results which we use in this article: first we cite two theorems from the theory of analytic combinatorics; then we establish an asymptotic expansion of the ratio c_n/c_{n-1} , where c_n is the number of connected chord diagrams.

In Section 3, we study the leading-log expansion and next-to^{*i*}-leading log expansions of the solution of (1). These are a way of organizing the double expansion of $G(x, L)$ and are relevant in quantum field theory. To do so, we are led to enumerate the connected chord diagrams C such that the first terminal chord is close to the last chord. We give recurrences that characterize these numbers and their exponential generating functions. We then establish asymptotic estimates on them. There will be implications on the physical level: we show that the dominant terms in all the log expansions only involve f_0 and f_1 , respectively the residue and the constant term of $F(\rho)$.

In Section 4, we study parameters on uniform large connected chord diagrams. We state a generic theorem that shows that numerous laws on connected chord diagrams obey to a Gaussian law, with a mean of the form $\lambda n + \mu \ln n$ and a logarithmic variance. In particular, we prove that the average number of terminal chords is $\ln n$. Finally, we see that the average index of the first terminal chord, a parameter which does not satisfy the above-mentioned theorem, is $2n/3$. The methods used are new and interesting as combinatorics.

2. ENUMERATIVE BACKGROUND

2.1. A BRIEF HISTORICAL BACKGROUND ON CONNECTED CHORD DIAGRAMS

Chord diagrams and their enumeration are not only relevant in quantum field theory; they also appear in various other areas of mathematics: knot theory [16, 3, 19] (in particular the Vassiliev invariants), graph sampling [1], analysis of computer structures [6], and even bioinformatics [10, 2].

Concerning more particularly the *connected* chord diagrams, Touchard seems to be the first person in 1952 to be interested in their enumeration [17]. More precisely, he characterized the number of connected diagrams with n chords and m crossings as a solution of a system of equations. Subsequently, Stein provided an explicit recurrence relation for the number of connected chord diagrams (but without considering this time the number of crossings) [14], as stated in the following proposition.

Proposition 5 (Stein [14]). *Let c_n be the number of connected diagrams with n chords. These numbers satisfy the relations $c_1 = 1$ and for $n \geq 2$,*

$$c_n = (n-1) \sum_{k=1}^{n-1} c_k c_{n-k}. \quad (3)$$

This formula has also been shown by Nijenhuis and Wilf, but this time thanks to a constructive combinatorial proof [13]. Let us also mention that (3) is equivalent to

$$c_n = \sum_{k=1}^{n-1} (2k-1) c_k c_{n-k}. \quad (4)$$

As for the asymptotic behaviour of the number c_n of connected diagrams with n chords, Stein and Everett gave the estimate

$$c_n \sim \frac{1}{e} (2n-1)!!$$

in [15]. In particular, since the number of (non necessarily connected) diagrams with n chords is $(2n-1)!!$, this implies that a large random chord diagram is connected with a probability e^{-1} .

Some decades later, Flajolet and Noy refined this result. Indeed, they proved in [7] that the number of connected components in a large random chord diagram (minus 1) follows a Poisson law of parameter 1. Moreover, they showed that if L_n denotes the size of the largest component in a random diagram with n chords, then $n - L_n$ is also distributed like a Poisson law of parameter 1.

Recently, Michael Borinsky computed in [4] an asymptotic expansion of the number of connected diagrams c_n , along with cardinals of similar objects.

2.2. PRELIMINARIES ON ANALYTIC COMBINATORICS

The majority of our proofs are based on the reference book by Flajolet and Sedgewick [8]. We present here the two main analytic combinatorics theorems of this paper.

First of all, let us mention that we use in this document two different notions of generating function. Given a sequence a_n , the *ordinary generating function* of the numbers a_n is defined as $\sum_{n \geq 0} a_n z^n$, while the *exponential generating function* is defined as $\sum_{n \geq 0} a_n z^n / n!$. Both notions have their advantages and drawbacks, especially when we try to enumerate chord diagrams. That is why we will juggle the two notions.

The first theorem, maybe the most representative of the theory, is called the *transfer theorem*. It relates the singular expansion of the series and the asymptotic behaviour of its coefficients.

Theorem 6 (Transfer Theorem). *Consider Δ a complex domain of the form*

$$\{z \mid |z| < \rho + \varepsilon\} \cap \{z \mid |\operatorname{Arg}(z - \rho)| > \alpha\},$$

with $\varepsilon > 0$ and $\alpha \in (0, \pi/2)$.

Let $f(z) = \sum_{n \geq 0} f_n z^n$ be an analytic function on Δ . If the singular behaviour of f in the vicinity of ρ is

$$f(z) \underset{z \rightarrow \rho}{\sim} c (1 - z/\rho)^{-\alpha} \ln(1 - z/\rho)^\beta,$$

where α is a complex number which does not belong to $\mathbb{Z}_{<0}$, β any integer and c a non zero constant, then

$$f_n \sim \frac{c}{\Gamma(\alpha)} \rho^{-n} n^{\alpha-1} (\ln n)^\beta.$$

In this article, the analyticity of our functions on a domain with the same shape as Δ is generally obvious (mainly because we have explicit expressions). The justification of analyticity will be then omitted, except if there is a subtlety to stress.

Example of use of transfer theorem. Consider the exponential generating function of $(2n-3)!!$. We will prove that it is equal to $1 - \sqrt{1-z}$. By the transfer theorem, the numbers $(2n-3)!!/n!$ are equivalent to $-n^{-3/2}/\Gamma(-1/2) = n^{-3/2}/(2\sqrt{\pi})$. This can be checked by the Sterling formula.

The next theorem deals with the Quasi-Powers Theorem, stated under a form which will be useful for us.

Theorem 7 (Theorem IX.11 of [8] – Quasi-Powers Theorem). *Let $F(z, u) = \sum_{n \geq n_0, k \geq 0} f_{n,k} z^n u^k$ be a bivariate function with non-negative coefficients such that $\sum_{k \geq 0} f_{n,k} = 1$ for $n \geq n_0$.*

(i) Analytic representation. *The function $F(z, u)$ admits the representation*

$$F(z, u) = A(z, u) + B(z, u)(1 - z)^{-\alpha(u)}$$

where $A(z, u)$ and $B(z, u)$ are analytic on a domain of the form $\{(z, u) \mid |z| \leq r, |u - 1| < \varepsilon\}$, with $r > 1$ and $\varepsilon > 0$. Assume also that $\alpha(u)$ is analytic at 1 such that $\alpha(1)$ is not a non-positive integer and $B(1, 1) \neq 0$.

(ii) Variability condition. *One has $\alpha'(1) + \alpha''(1) \neq 0$.*

Then the random variable X_n such that $\mathbb{P}(X_n = k) = f_{n,k}$ converges in distribution to a Gaussian variable. The corresponding mean is $\alpha'(1) \ln n$ and the variance is $(\alpha'(1) + \alpha''(1)) \ln n$.

2.3. A REFINED ASYMPTOTIC RESULT

We will need a precise asymptotic expansion of c_{n-1}/c_n . The dominant term has already been given by Stein and Everett: they showed that $c_{n-1}/c_n \sim 1/2n$. This expansion can also be deduced from the article of Michael Borinsky [4].

Proposition 8. *The ratio between the numbers of connected chord diagrams with $n - 1$ arcs and n arcs is asymptotically equivalent to*

$$\frac{c_{n-1}}{c_n} = \frac{1}{2n} + \frac{1}{4n^2} - \frac{1}{2n^3} - \frac{29}{8n^4} + O\left(\frac{1}{n^5}\right).$$

The proof simply relies on what is often called *bootstrapping*. We need first to establish a lemma which bounds the contribution of the central terms in the sum $(n - 1) \sum_{k=1}^{n-1} \frac{c_k c_{n-k}}{c_n}$.

Lemma 9. *If c_n denotes the number of connected diagrams with n chords, we have for fixed $j \geq 2$ the estimate*

$$(n - 1) \sum_{k=j}^{n-j} \frac{c_k c_{n-k}}{c_n} = O\left(\frac{1}{n^{j-1}}\right).$$

Proof. The core of the proof lies in the inequality

$$(2n - 1) c_{n-1} < c_n < 2n c_{n-1}, \quad (5)$$

holding for all $n \geq 5$. This has been stated by Stein and Everett in [15, Lemmas 3.1 and 3.4] and was proved by a (technical) induction. Notice that this inequality justifies the previous estimate $c_{n-1}/c_n \sim 1/2n$. For $n \geq 8$ and $k \in \{5, \dots, \lfloor \frac{n}{2} \rfloor\}$, we then have

$$c_k c_{n-k} \leq \frac{2k}{2n - 2k + 1} c_{k-1} c_{n-k+1} \leq c_{k-1} c_{n-k+1}.$$

The inequality $c_k c_{n-k} \leq c_{k-1} c_{n-k+1}$ is also true for $k \in \{2, 3, 4\}$ since $c_k/c_{k-1} \leq c_4/c_3 < 7$ and thanks to (5), we have $c_{n-k+1}/c_{n-k} \geq (2n - 2k + 1) \geq 9$ for every $n \geq 8$ and $k \in \{2, 3, 4\}$ (Inequality (5) applies because $n - k + 1 \geq 5$). Therefore, we show by a basic induction that for every $n \geq 8$ and for $2 \leq k \leq j \leq \lfloor \frac{n}{2} \rfloor$,

$$c_k c_{n-k} \leq c_{k-1} c_{n-k+1} \leq \dots \leq c_{j-1} c_{n-j+1} \leq c_j c_{n-j}. \quad (6)$$

Taking in particular the inequality $c_k c_{n-k} \leq c_{j-1} c_{n-j+1}$ and summing over $k \in \{j+1, n-j+1\}$, we obtain

$$\sum_{k=j+1}^{n-j-1} c_k c_{n-k} \leq n c_{j+1} c_{n-j-1}.$$

Consequently,

$$(n-1) \sum_{k=j}^{n-j} \frac{c_k c_{n-k}}{c_n} = O\left(n \frac{c_{n-j}}{c_n} + n^2 \frac{c_{n-j-1}}{c_n}\right).$$

But c_{n-j}/c_n can be written as $\prod_{i=1}^j c_{n-i}/c_{n-i+1}$, so is equivalent to $1/(2n)^j$ (we have used the estimate $c_{n-1}/c_n \sim 1/2n$). Plugging this in the previous equality directly gives the lemma. \square

Now let us describe how to find an expansion of c_{n-1}/c_n .

Proof of Proposition 8. By combining the previous lemma and (3), we deduce that

$$1 = 2(n-1) \sum_{k=1}^j \frac{c_k c_{n-k}}{c_n} + O\left(\frac{1}{n^j}\right)$$

holds for $j \geq 1$. For $j = 1$, we recover $c_{n-1}/c_n \sim 1/2n$. For $j = 2$, we have

$$1 = 2(n-1) \frac{c_{n-1}}{c_n} + 2(n-1) \frac{c_{n-2} c_{n-1}}{c_{n-1} c_n} + O\left(\frac{1}{n^2}\right).$$

Setting $c_{n-1}/c_n = 1/2n + d_n$ leads to

$$0 = -\frac{1}{n} + 2(n-1)d_n + 2(n-1) \left(\frac{1}{2(n-1)} + o\left(\frac{1}{n}\right) \right) \left(\frac{1}{2n} + o\left(\frac{1}{n}\right) \right) + o\left(\frac{1}{n}\right)$$

(we have used the fact that $d_n = o(1/n)$), and so $d_n \sim 1/4n^2$.

This process can be repeated for $j = 3, 4, \dots$ to find the predicted expansion of c_{n-1}/c_n . \square

3. LOG EXPANSIONS OF THE SOLUTION OF THE DYSON-SCHWINGER EQUATION

3.1. CONTEXT

In this section, we show how we can deduce from Theorem 4 asymptotic properties on the *log expansions* in quantum field theory. Let us explain first what is a log expansion.

Suppose we have an expansion with the following form

$$G(x, L) = 1 + \sum_{i \geq 1} \sum_{j \geq i} a_{i,j} L^i x^j.$$

The particular Dyson-Schwinger equations we are interested in have their solution in this form as do a broad class of perturbative expansions in quantum field theory. Then, rather than thinking of the sum first as an expansion in one of the variables with coefficients which are series in the other variables, we can take an expansion which takes variables together.

Specifically, we can write the expansion as

$$G(x, L) = \sum_{k \geq 0} \sum_{i \geq 0} a_{i,i+k} (Lx)^i x^k.$$

The $k = 0$ part of this sum, namely the terms of $G(x, L)$ where the powers of L and x are the same, is known as the *leading log expansion*, the $k = 1$ part of this sum, namely the terms of $G(x, L)$ where the power of x is one more than the power of L is known as the *next-to-leading log expansion*. The $k = 2$ is known as the *next-to-next-to-leading log expansion* and so on.

This leading log language comes from the fact that L is the logarithm of some appropriate energy scale, while x is the coupling constant which is treated as a small parameter. So the leading log expansion captures the maximal powers of x relative to the powers of the energy scale, and so is in an important sense the leading term. The next-to-leading log expansion is the next part; it is suppressed by one power of x , and so on.

Furthermore the full log expansion is algebraically and analytically meaningful in the sense that the contributions of larger primitive graphs and new (presumably) transcendental numbers appear further out in the next-to-next-to... hierarchy. We see this manifested in our results, but it is a much more general physical fact (compare [11]).

In view of (2), the leading log expansion for the Dyson-Schwinger equations is

$$- \sum_{\substack{C \\ b(C)=|C|}} \frac{(-Lx)^{|C|}}{|C|!} f_0^{|C|}$$

while the next-to-leading log expansion is

$$\begin{aligned} - \sum_{\substack{C \\ b(C)=|C|-1}} \frac{(-Lx)^{|C|-1}x}{(|C|-1)!} f_0 f_0^{|C|-2} f_1 - \sum_{\substack{C \\ b(C)=|C|}} \frac{(-Lx)^{|C|-1}x}{(|C|-1)!} f_1 f_0^{|C|-1} \\ = - \sum_{\substack{C \\ b(C) \geq |C|-1}} \frac{(-Lx)^{|C|-1}x}{(|C|-1)!} f_0^{|C|-1} f_1 \quad (7) \end{aligned}$$

In general the next-to- i -leading log expansion is

$$- \sum_{\substack{C \\ b(C) \geq |C|-i}} \frac{(-Lx)^{|C|-i}x^i}{(|C|-i)!} f_{b(C)-|C|+i} f_0^{|C|-k} \prod_{j=2}^k f_{t_j-t_{j-1}}$$

where $b(C) = t_1 < t_2 < \dots < t_k$ are the terminal chords of C . We switch the signs from now on, both overall and of L , because they are the result of the conventions of [18] and perhaps not actually a good choice.

All this suggests that it is worthwhile to study connected diagrams C such that $b(C) \geq |C|-i$, where i is fixed. The present section continues by establishing numerous enumerative and asymptotic results concerning these diagrams.

3.2. RECURRENCE EQUATIONS

We begin by an induction that characterizes the number of connected diagrams C of size n such that every terminal chord has index between $n-k$ and n .

Proposition 10. *Fix $k \geq 0$. For $n > 0$, let $b_{n,k}$ be the number of connected diagrams C with n chords such that $b(C) \geq n-k$, and $c_n = b_{n,n-1}$ the number of connected diagrams with n chords. For every $k \geq 0$, we have $b_{1,k} = 1$ and for $n \geq 2$*

$$b_{n,k} = (2n-3)b_{n-1,k} + \sum_{i=1}^{\min(k,n-2)} (2i-1)c_i b_{n-i,k-i}. \quad (8)$$

This recurrence relation enables to compute the first values of $b_{n,k}$:

$$\begin{aligned} b_{1,0} &= 1, b_{2,0} = 1, b_{3,0} = 3, b_{4,0} = 15, b_{5,0} = 105, b_{6,0} = 945, \dots \\ b_{1,1} &= 1, b_{2,1} = 1, b_{3,1} = 4, b_{4,1} = 23, b_{5,1} = 176, b_{6,1} = 1689, \dots \\ b_{2,1} &= 1, b_{2,2} = 1, b_{3,2} = 4, b_{4,2} = 27, b_{5,2} = 221, b_{6,2} = 2210, \dots \end{aligned}$$

Proof. Equation (8) is derived from a specific decomposition of the connected chord diagrams, which we are going to describe, and which is illustrated by Figure 3. This recursion has a good transcription in terms of exponential generating functions (see Proposition 12).

Let C be a connected chord diagram of size n such that $b(C) \geq n-k$. If we remove the root chord of C , two exclusive possibilities can occur:

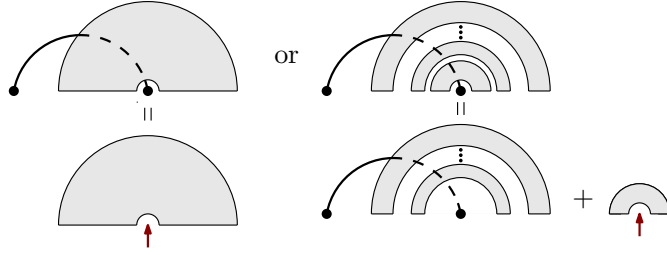


FIGURE 3. Illustration of the decomposition from the proof of Proposition 10.

- **The obtained diagram C' is still connected.** In this case, C' has $n - 1$ chords and we have $b(C') = b(C) - 1$, and so $b(C') \geq n - 1 - k$. Moreover, to recover the diagram C , we need (and it is sufficient) to remember the position of the right endpoint of the root chord of C . Since C' has $n - 1$ chords, there are $2(n - 1) - 1$ possible positions. That is why the number of such diagrams C is given by $(2n - 3)b_{n-1,k}$.
- **We obtain several connected components C_1, \dots, C_s .** We denote by i the number of chords in C_s , and by C' the (connected) diagram obtained by removing C_s from C . Diagram C' has $n - i$ chords, and the position of the first terminal chord has remained unchanged, hence $b(C') \geq (n - i) - (k - i)$. Observe then that it is possible to recover C from C' and C_s provided the position of the root chord through C_s (there are $(2i - 1)$ such possible positions). The number of such diagrams C is thus $(2i - 1)c_i b_{n-i,k-i}$.

The conjunction of these two cases infers Equation (8). \square

Remark that for $k = 0$, Recurrence (8) is simply $b_{n,0} = (2n - 3)b_{n-1,0}$. This provides a nice formula (and a combinatorial proof!) for the numbers $b_{n,0}$, which correspond to the numbers of connected chord diagrams with exactly one terminal chord (the last chord of a diagram is necessarily terminal).

Corollary 11. *The number of connected diagrams with n chords and only one terminal chord is $(2n - 3)!!$.*

The recurrence relation of Proposition 10 can be transformed in an effective way to compute the exponential generating functions of the numbers $b_{n,k}$:

Proposition 12. *Let $B_k(z)$ be the exponential generating function of the connected chord diagrams C such that $b(C) \geq n - k$. For every pair of integers i, k , we consider an i th antiderivative¹ $B_k^{[i]}$ for B_k , that is, a function $B_k^{[i]}$ such that its i th derivative is equal to B_k . There exists a constant β_k and a polynomial P_k of degree k such that*

$$B_k(z) = \sqrt{1 - 2z} \left(\beta_k + \sum_{i=1}^k (2i - 1) c_i \int_0^z (1 - 2x)^{-3/2} B_{k-i}^{[i-1]}(x) dx \right) + P_k(z) \quad (9)$$

where c_i is the number of connected diagrams with i chords.

Proof. This proof can be divided into two steps. First, we translate (8) in terms of the functions B_{k-i} with $i \in \{0, \dots, k\}$, which gives a first order differential equation in B_k . Then, we simply solve this differential equation.

¹It is possible to define uniquely this antiderivative by setting for example $\frac{\partial^j B_k^{[i]}}{\partial z^j}(0) = 0$ for $j \in \{0, \dots, i - 1\}$, but in practice, it is more convenient to take any antiderivative we find.

Let $n > k + 1$. Dividing (8) by $(n - 1)!$ and writing $2n - 3$ as $2(n - 1) - 1$ induces that

$$\frac{b_{n,k}}{(n-1)!} - 2 \frac{b_{n-1,k}}{(n-2)!} + \frac{b_{n-1,k}}{(n-1)!} = \sum_{i=1}^k (2i-1) c_i \frac{b_{n-i,k-i}}{(n-1)!}. \quad (10)$$

Observe in all generality that the series of general term $a_{n-i-1} z^{n-1}/(n-1)!$ is the $(-i)$ th derivative of the exponential generating function of the sequence a_n if i is non-positive, and an $(i-1)$ th antiderivative if i is positive. We then recognize in (10) the coefficients of z^{n-1} in $\frac{\partial B_k}{\partial z}, z \frac{\partial B_k}{\partial z}, B_k, B_{k-i}^{[i-1]}$, respectively. Thus (10) can be translated by

$$(1-2z) \frac{\partial B_k}{\partial z} + B_k = \sum_{i=1}^k (2i-1) c_i B_{k-i}^{[i-1]} + Q_k(z),$$

where $Q_k(z)$ is a polynomial of degree k whose presence is due to the fact that the first coefficients of $B_{k-i}^{[i-1]}$ can vary, but also because (10) holds only for $n > k + 1$.

We can solve this differential equation quite straightforwardly: we divide by $(1-2z)^{3/2}$ both sides and recognize from the left side the derivative of $(1-2z)^{-1/2} B_k$. Integrating this equation then leads to (9). (We have set $P_k(z) = \sqrt{1-2z} \int_0^z (1-2x)^{-3/2} Q_k(x) dx$. Some easy calculus shows that P_k is also a polynomial of degree k .) \square

Remark 1. It is simple to compute the series $B_k(z)$ by recursion thanks to Formula (9). The method is the following: for each i , we begin by compute the antiderivatives $B_{k-i}^{[i-1]}$, then plug them into (9), evaluate the formula and then eliminate β_k and $P_k(z)$ thanks to the first values of $B_k(z)$ given by Proposition 10. We thus obtain:

$$\begin{aligned} B_0(z) &= 1 - \sqrt{1-2z}, \\ B_1(z) &= 1 + z + \frac{1}{2} \sqrt{1-2z} \ln(1-2z) - \sqrt{1-2z}, \\ B_2(z) &= \left(\frac{\ln(1-2z)}{2} - \frac{\ln(1-2z)^2}{8} + z - 3 \right) \sqrt{1-2z} + 3 - 2z + \frac{z^2}{2}. \end{aligned}$$

It is important to notice that the method is automatic. In this regard, a `maple` file is available along with the `arXiv` version of this paper.

Remark 2. The foregoing gives information about the generating function of connected diagrams C such that $b(C) \geq n-k$ but nothing about the distribution of $f_{b(C)-i} f_0^{|C|-k} \prod_{j=2}^k f_{t_j-t_{j-1}}$ in the leading-log coefficients. However it is easy to adapt the same approach to enumerate diagrams to specific cases where the $t_j - t_{j-1}$ are fixed.

Example. Let us consider $A(z)$ the exponential generating function of connected diagrams such that the only terminal chords are the third to last and last ones (i.e. connected diagrams C such that $b(C) = t_1 = |C| - 2$ and $t_2 = |C|$), and let us use the same decomposition as in the proof of Proposition 10. Removing the root chord in such diagrams leads to two possibilities:

- The resulting diagram has only one component; starting from the end, the positions of the terminal chords do not change.
- It has several components. Since each component necessarily has at least one terminal chord, the number of components is exactly two: the top component has only one terminal chord, and the bottom component is the only connected diagram with 2 chords (with 3 possibilities of insertion for the root chord).

This consideration leads to the recurrence

$$a_n = a_{n-1} + 3(2n-7)!!,$$

where a_n is the number of connected diagrams with n chords such that the only terminal chords are the third to last and last ones. By the same process than previously, we can then prove that

$$A(z) = (z-1)\sqrt{1-2z} + \frac{z^2}{2} - 2z + 1.$$

In all generality, similar recursions exist for diagrams where the gaps between the terminal chords $t_2 - t_1, \dots, t_k - t_{k-1}$ are given, but equations are more tedious to state (although the method will fundamentally remain the same). If the reader would like to compute such generating functions, a procedure is written in the aforementioned `maple` file.

3.3. ASYMPTOTIC BEHAVIOUR

Now that we have stated how to compute the numbers $b_{n,k}$, we are interested by their asymptotic behaviour. The following theorem gives the asymptotic estimate.

Theorem 13. *The number of connected diagrams C with n chords such that $b(C) \geq n - k$ is asymptotically equivalent to*

$$\frac{1}{\sqrt{\pi} 2^{k+1} k!} \frac{\ln(n)^k 2^n n!}{n^{3/2}}.$$

Proof. We just apply the transfer theorem (Theorem 6) to the exponential generating function B_k characterized by the following lemma. Indeed this lemma shows that $B_k(z) - B_k(1/2)$ is equivalent to $\frac{(-1)^{k+1}}{2^k k!} \sqrt{1-2z} \ln(1-2z)^k$ when $z \rightarrow 1/2$. \square

Lemma 14. *The exponential generating function B_k from Proposition 12 is a polynomial in terms of $\sqrt{1-2z}$ and $\ln(1-2z)$:*

$$B_k(z) = B_k\left(\frac{1}{2}\right) + \sqrt{1-2z} Q_k\left(\sqrt{1-2z}, \ln(1-2z)\right), \quad (11)$$

where $Q_k(x, y)$ is a polynomial of degree k in y such that the coefficient of y^k is

$$[y^k] Q_k(x, y) = \frac{(-1)^{k+1}}{2^k k!}$$

(there is no term in x^i , with $i \geq 1$).

Proof. This lemma can be shown by induction on k .

For $k = 0$, we have seen that $B_0(z) = 1 - \sqrt{1-2z}$.

For $k \geq 1$, we have to check that the statement of the lemma is compatible with (9). First, using the induction hypothesis, we verify that any $(i-1)$ th antiderivative $B_{k-i}^{[i-1]}$ of B_i for $i \in \{1, \dots, k\}$ is a polynomial in $\sqrt{1-2z}$ and $\ln(1-2z)$ such that the degree in $\ln(1-2z)$ does not exceed $k-i$. (To compute an antiderivative of $(1-2z)^{i/2} \ln(1-2z)^j$, we repeatedly integrate by parts using the equality

$$(i+2) \int_0^z (1-2x)^{i/2} \ln(1-2x)^j dx = -(1-2z)^{i/2+1} \ln(1-2z)^j - 2j \int_0^z (1-2x)^{i/2} \ln(1-2x)^{j-1} dx$$

until the degree in $\ln(1-2x)$ in the integrand reaches 0.) Using (9), it is then not hard to check that $B_k(z)$ can be put into the form (11). If we search in (9) for what could contribute to the term in y^k in $Q_k(x, y)$, we realize that the only possibility comes from the monomial $\sqrt{1-2z} \ln(1-2z)^{k-1}$ in $B_{k-1}(z)$. Indeed, we can observe that

$$\sqrt{1-2z} \int_0^z (1-2x)^{-3/2} \times \sqrt{1-2x} \ln(1-2x)^{k-1} dx = \frac{-1}{2k} \sqrt{1-2z} \ln(1-2z)^k.$$

By recurrence, we know that the coefficient of $\sqrt{1-2z} \ln(1-2z)^{k-1}$ in $B_{k-1}(z)$ is $\frac{(-1)^k}{2^{k-1}(k-1)!}$, so the coefficient of $\sqrt{1-2z} \ln(1-2z)^k$ in $B_{k-1}(z)$ must be $\frac{-1}{2k} \times \frac{(-1)^k}{2^{k-1}(k-1)!} = \frac{(-1)^{k+1}}{2^k k!}$. \square

Once again, the foregoing does not give any information about the asymptotic distribution on the terminal chords. However we can recover it by repeating the same reasoning for the number of connected diagrams such that only the last k chords for the intersection order are terminal, and observe that the asymptotic behaviour is identical.

Theorem 15. *The number $o_{n,k}$ of connected diagrams C with n chords such that the only terminal chords are the last k chords is asymptotically equivalent to*

$$\frac{1}{\sqrt{\pi} 2^{k+1} k!} \frac{\ln(n)^k 2^n n!}{n^{3/2}}.$$

Proof. (Sketch.) Using the decomposition of the proof of Proposition 10, we find that the numbers $o_{n,k}$ satisfy

$$o_{n,k} = (2n-3) o_{n-1,k} + o_{n-1,k-1}.$$

This recurrence relation can be then translated into the differential equation

$$(1-2z) \frac{\partial O_k}{\partial z} + O_k = O_{k-1} + \tilde{Q}_k(z),$$

where $O_k(z) = \sum_{n \geq 0} o_{n,k} z^n / n!$ and \tilde{Q}_k a polynomial of degree k . Its solutions can be put into the form

$$O_k(z) = \sqrt{1-2z} \left(\tilde{\beta}_k + \int_0^z (1-2x)^{-3/2} O_{k-1}(x) dx \right) + \tilde{P}_k(z),$$

where $\tilde{\beta}_k$ is a constant and $\tilde{P}_k(z)$ a polynomial. By recurrence, we can then prove that O_k is a polynomial in $\sqrt{1-2z}$ and $\ln(1-2z)$ such that the contributing term for the singularity analysis is $\frac{(-1)^{k+1}}{2^k k!} \sqrt{1-2z} \ln(1-2z)^k$. We recover the expected asymptotic regime by the transfer theorem. \square

The consequence of the similarity between Theorem 13 and Theorem 15 will be described in the next subsection.

3.4. APPLICATION TO THE LOG EXPANSIONS

The leading log expansion is particularly simple because it only counts chord diagrams where only the last chord is terminal. By Corollary 11, these are easy to count, and the monomial in the f_i is simply a power of f_0 . Therefore it suffices to understand $B_0(z)$. Specifically, the leading log expansion is

$$B_0(Lx f_0) = 1 - \sqrt{1-2Lx f_0} \quad (12)$$

The next-to-leading log expansion is not too difficult either. By (7) it suffices to understand $B_1(z)$. Note, however, that the power of $Lx f_0$ is $|C|-1$, and we are dividing by $(|C|-1)!$, so the next-to-leading log expansion is actually given in terms of the derivative of $B_1(z)$. Specifically, the next-to-leading log expansion is

$$\frac{d}{dz} B_1(z)|_{z=Lx f_0} x f_1 = x f_1 \left(1 + \frac{1}{\sqrt{1-2Lx f_0}} \ln \left(\frac{1}{\sqrt{1-2Lx f_0}} \right) \right) \quad (13)$$

The next-to-next-to-leading log expansion is a bit more complicated. Here we are considering any chord diagram with $b(C) \geq |C|-2$. Now there are different possible monomials. If all of the last three chords are terminal then we get $f_0 f_0^{|C|-3} f_1^2$ while if only the last and the third last are terminal we get $f_0 f_0^{|C|-2} f_2$. If only the last two chords are terminal we get $f_1 f_0^{|C|-2} f_1$

and if only the last chord is terminal we get $f_2 f_0^{|C|-1}$. All together two different monomials appear, $f_2 f_0^{|C|-1}$ in the case that either the last and third last or just the last are terminal, and $f_1^2 f_0^{|C|-2}$ in the case that either the last two or the last three are all terminal. In all cases we will need to take two derivatives since the powers and factorials are in terms of $|C| - 2$ for the next-to-next-to leading log expansion rather than in terms of $|C|$ for the exponential generating functions B_k .

Using the $A(z)$ from the example in Subsection 3.2 we can calculate the next-to-next-to-leading log expansion explicitly:

$$\begin{aligned} & x^2 f_2 f_0 \left(\frac{d^2}{dx^2} (A(z) + B_0(z)) \right) \Big|_{z=Lx f_0} + x^2 f_1^2 \left(\frac{d^2}{dx^2} (B_2(z) - A(z) - B_0(z)) \right) \Big|_{z=Lx f_0} \\ &= x^2 f_0 f_2 \left(1 + \frac{3Lx f_0}{(1 - 2Lx f_0)^{3/2}} \right) + \frac{x^2 f_1^2 (\ln(1 - 2Lx f_0) - 4) \ln(1 - 2Lx f_0)}{8(1 - 2Lx f_0)^{3/2}}. \end{aligned}$$

Latter log expansions work similarly.

Let us compare these results to the results of Krüger and Kreimer in [11]. Their methods are also combinatorial but are quite different. They are based on words on the alphabet of primitive graphs operated on by shuffle and Lie bracket. Despite these differences we are both modelling the same underlying physics, so our answers should agree on the common domain of applicability.

Our results correspond to their Yukawa case with only one primitive. In fact we deal with any Dyson-Schwinger equation with this shape. They could also do so, but chose to only make the Yukawa and QED examples explicit. On the other hand their work is more general in that they deal with any number of primitives in the Yukawa and QED example. In view of [9] our results should also generalize to any number of primitives and to Dyson-Schwinger equations of QED shape and other shapes (corresponding to different s parameters in the setup of [18]). This will be worked out in the future.

Our leading log calculations are, as they should be, identical (compare (12) to Equation 221 of [11]). The next-to-leading log (compare (13) and Equation 227 of [11]) are very similar. First in our case we are not considering a new primitive graph at 2 loops, which in the language of Krüger and Kreimer would say that $\Phi_R(\Gamma_2) = 0$. The second thing to notice is a spurious 1 in the derivative of $B_1(z)$. Its presence is due to different boundary conditions. They explicitly set their generating function to have no constant term (see the line after Equation 146) while our boundary conditions are determined by the chord diagrams: this particular 1 corresponds to the connected chord diagram with two chords. Finally, note that they have a more complicated expression in place of our f_1 . In both cases this number is the new period. Krüger and Kreimer call it $\Theta(a_1, a_1)$; they note that it cannot be canonically identified with a single Feynman graph. From our perspective we see it naturally as the next term in the expansion for the original primitive.

Turning to the next-to-next-to-leading log expansion, we again see that our solution is built on of the same kinds of pieces as theirs. Their greater generality shows up more strongly here as our solution is strictly simpler. We also see more clearly at this level how our different perspectives result in different characterizations of the new primitives.

What are the benefits and disadvantages of our techniques compared to the techniques of Krüger and Kreimer? Both methods have a combinatorially derived master equation that determines everything. For us this would be the recursive decomposition of the previous sections – we did not write it out as an integral or differential equation in general (only in the important special case of Proposition 12), but the example in Subsection 3.2 illustrates how it works in general. For both groups the master equation is not fully explicit. In Krüger and Kreimer's set up this manifests itself in the dependence on matrix bracket coefficients for which it is unclear

how automatically or rapidly they can be computed. The lack of explicitness however has a different flavour in each case coming from the different combinatorial objects.

Our technique also differs in how it indexes the periods which contribute to the expansions. Krüger and Kreimer tie them to individual graphs where possible and treat the others, coming from their Θ -expressions on the same level. We do not give these periods individual meanings but see them as coming from later terms in the expansion our one primitive; this makes our periods less combinatorial, but they are organized into tidy monomials so one can better see the different pieces that build them. Here both techniques have advantages and one would hope to play them off each other to get an even better understanding. The same can be said about the different underlying combinatorial frameworks – the physics is described both by our chord diagrams and by their words and it is not obvious, but is potentially useful, that both these objects describe the same underlying structures.

Finally, we can consider the significance of the asymptotic results of Subsection 3.3 to the log expansions. Here something very interesting happens. The chord diagrams where the only terminal chords are the last k chords dominate completely in the sense that as $n \rightarrow \infty$ almost all chord diagrams with n chords and $b(C) \geq n - k$ have the last k chords terminal. What this means is that provided $F(\rho)$ is not outrageous (eg the f_i are bounded) we should expect the chord diagrams with the last k chords terminal to completely determine the asymptotic behaviour of the next-to- k -leading log expansion. That is, the next-to- k -leading log expansion should behave as if all chord diagrams contribute the monomial $f_0^{|C|-k+1} f_1^{k-1}$, so the asymptotic behaviour of the next-to- k -leading log expansion is given by

$$\frac{d^k}{dz^k} B_k(z)|_{z=Lx f_0 f_1^{k-1}} \quad (14)$$

This is nice for two reasons. First it says that the other f_i are not playing a significant role asymptotically – this means only two numbers, f_0 and f_1 , are controlling their asymptotic behaviours. Second the master equation to generate the $B_k(z)$ is fairly simple and can be computed fully automatically. This is much simpler than the situation for chord diagrams with specific gap patterns as we calculated for the next-to-next-to-leading log expansion and contains no mysteries which require human intervention to compute.

4. STATISTICS ON TERMINAL CHORDS

4.1. STATEMENT OF THE META-THEOREM AND EXAMPLES

In this section, we study several statistics concerning terminal chords in connected diagrams, such like its numbers, the number of terminal chords that are consecutive for the intersection order, etc. We establish a meta-theorem that shows that a lot of random variables on connected chord diagrams have a Gaussian limit law with logarithmic variance.

Before stating this theorem, we need to define three subsets of connected chords diagrams based on the shape of the diagram obtained by removing the root chord:

- $\mathcal{C}_1 = \{C \mid \text{Removing the root chord from } C \text{ leads to a unique component } C_1.\}$,
- $\mathcal{C}_2 = \{C \mid \text{Removing the root chord leads to a chord and a component } C_2, \text{ top to bottom.}\}$,
- $\mathcal{C}_3 = \{C \mid \text{Removing the root chord leads to a component } C_3 \text{ and a chord, top to bottom.}\}$.

These three subsets are illustrated by Figure 4.

Theorem 16. *Let $\lambda_1, \lambda_2, \lambda_3$ be three integers (not all equal). Consider a function f on connected chord diagrams such that for every $C \in \mathcal{C}_i$ with $i = 1, 2, 3$, $f(C) = f(C_i) + \lambda_i$ (see above for the definition of $\mathcal{C}_1, \mathcal{C}_2, \mathcal{C}_3$). If X_n denotes a random chord diagram of size n under the uniform distribution, then*

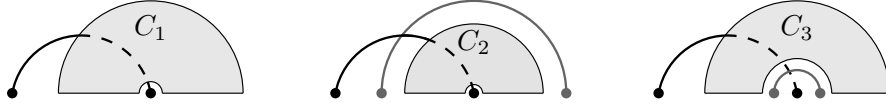


FIGURE 4. From left to right, a schematic representation of an element of \mathcal{C}_1 , \mathcal{C}_2 , \mathcal{C}_3 .

$f(X_n)$ is a random variable such that $\frac{f(X_n) - \lambda_1 n - \mu \ln(n)}{\sigma \sqrt{\ln(n)}}$ converges in distribution to a standard Gaussian law, where

$$\mu = \frac{\lambda_2}{2} + \frac{\lambda_3}{2} - \lambda_1, \quad \sigma^2 = \frac{(\lambda_2 - \lambda_1)^2}{2} + \frac{(\lambda_3 - \lambda_1)^2}{2}.$$

Among other things, this theorem implies that the relevant diagrams under the uniform distribution are those whose recursive decomposition only uses diagrams from \mathcal{C}_1 , \mathcal{C}_2 and \mathcal{C}_3 . The other diagrams are asymptotically negligible (the proof of this theorem just uses this fact).

Let us illustrate Theorem 16 with some examples. If we denote by T_n the random variable on connected diagrams with n chords that counts the terminal chords, we can see that $T_n = f(X_n)$ where f and X_n are described in the statement of the theorem with $\lambda_1 = 0$ and $\lambda_2 = \lambda_3 = 1$. (Only diagrams from \mathcal{C}_2 and \mathcal{C}_3 have a decomposition which induces terminal chords – the terminal chords correspond to the dark-grey ones in Figure 4.) Consequently, we have the following corollary.

Corollary 17. *The random variable T_n for the number of terminal chords asymptotically has a Gaussian limit law with a mean and a variance equivalent to $\ln(n)$.*

Now let us consider $G_{1,n}$, the random variable on connected diagrams with n chords that counts the pairs of terminal chords that are adjacent in the intersection order. Equivalently, $G_{1,n}$ counts the number of terminal chords c such that the chord that precedes c in the intersection order is also terminal. We can then notice that decompositions of diagrams from \mathcal{C}_1 and \mathcal{C}_2 do not induce such terminal chords (for the former, the only apparent chord is not terminal; for the latter, the chord that precedes the terminal chord is the root chord, which is not terminal), while decompositions for \mathcal{C}_3 do (the chord that precedes the terminal chord is the last chord of \mathcal{C}_3 which is terminal – the last chord of a connected diagram is always terminal). Therefore, we have $G_{1,n} = f(X_n)$ with $\lambda_1 = \lambda_2 = 0$ and $\lambda_3 = 1$, which gives the following result.

Corollary 18. *The random variable $G_{1,n}$ for the pairs of terminal chords that are adjacent for the intersection order has a Gaussian limit law with a mean and a variance asymptotically equivalent to $\frac{\ln(n)}{2}$.*

Remark. It is worth noting that the standard theory cannot be used directly. Indeed, the main obstacle is the non-analyticity of the ordinary generating functions, which *a priori* prevents any use of complex analysis. For instance, if we consider $C(z, u)$ the generating function of connected diagrams, where z refers to the number of chords and u to the number of terminal chords, this series satisfies the differential equation (which can be established by a straightforward combinatorial specification – see [8])

$$C(z, u) = zu + z \frac{2z \frac{\partial C}{\partial z}(z, u) - C(z, u)}{1 - 2z \frac{\partial C}{\partial z}(z, u) + C(z, u)}.$$

We can solve this non-linear differential equation (to some extent – the solution can be implicitly defined in terms of the Whittaker functions) but it seems to be impossible to deduce anything from there.

4.2. PROOF OF THEOREM 16

First of all, remark that we can assume that $\lambda_1 = 0$ without any loss of generality. Indeed, we can study $\tilde{f}(X_n) := f(X_n) - \lambda_1 n$ instead of $f(X_n)$. The new function \tilde{f} satisfies the conditions of Theorem 16 where the new set of parameters $\tilde{\lambda}_1, \tilde{\lambda}_2, \tilde{\lambda}_3$ is equal to 0, $\lambda_2 - \lambda_1, \lambda_3 - \lambda_1$. From the rest of this subsection, we suppose $\lambda_1 = 0$.

Before presenting the idea of the proof of Theorem 16, let us state an asymptotic equation governing the probabilities $\mathbb{P}(f(X_n) = k)$, that we denote shorthand $p_{n,k}$.

Lemma 19. *Let f and X_n as stated by Theorem 16 with $\lambda_1 = 0$, and let $p_{n,k}$ denote $\mathbb{P}(f(X_n) = k)$. Then, when n goes to infinity,*

$$p_{n,k} = (1 - n^{-1}) p_{n-1,k} + \frac{n^{-1}}{2} (p_{n-2,k-\lambda_2} + p_{n-2,k-\lambda_3}) + O(n^{-2}). \quad (15)$$

Proof. Under the condition $X_n \in \mathcal{C}_1$, the probability that $f(X_n)$ equals k is $p_{n-1,k}$. Indeed, removing the root chord from X_n gives a *uniform* connected diagram C_1 with $n-1$ chords such that $f(X_n) = f(C_1)$ (since $\lambda_1 = 0$). Similarly, the probability that $f(X_n)$ equals k under the condition $X_n \in \mathcal{C}_i$, with $i = 2$ or 3 , is $p_{n-2,k-\lambda_i}$. This shows that

$$\begin{aligned} p_{n,k} &= \mathbb{P}(X_n \in \mathcal{C}_1) p_{n-1,k} + \mathbb{P}(X_n \in \mathcal{C}_2) p_{n-2,k-\lambda_2} \\ &\quad + \mathbb{P}(X_n \in \mathcal{C}_3) p_{n-2,k-\lambda_3} + \mathbb{P}\left(X_n \notin \bigcup_{i=1,2,3} \mathcal{C}_i \text{ and } f(X_n) = k\right) \end{aligned}$$

We can then see that $\mathbb{P}(X_n \in \mathcal{C}_1) = (2n-3) \frac{c_{n-1}}{c_n}$ since the number of diagrams of size n in \mathcal{C}_1 is equal to the number of size of connected diagram of size $n-1$ (namely c_{n-1}) times the number of ways of inserting the root chord in this diagram ($2n-3$ ways to do it). By Proposition 8, we deduce that $\mathbb{P}(X_n \in \mathcal{C}_1) = 1 - n^{-1} + O(n^{-2})$. Similarly, $\mathbb{P}(X_n \in \mathcal{C}_2) = \mathbb{P}(X_n \in \mathcal{C}_3) = (2n-5) \frac{c_{n-2}}{c_n} = n^{-1}/2 + O(n^{-2})$. Finally, we deduce

$$\mathbb{P}\left(X_n \notin \bigcup_{i=1,2,3} \mathcal{C}_i \text{ and } f(X_n) = k\right) \leq 1 - \sum_{i=1,2,3} \mathbb{P}(X_n \in \mathcal{C}_i) = O(n^{-2}),$$

which proves the lemma. \square

Lemma 19 suggests that the recursive equation relating the numbers $p_{n,k}$ is easy to study (mainly because it *almost* involves polynomial coefficients), but the presence of the error term $O(n^{-2})$ makes the analysis tricky. The idea then consists in forgetting this term and studying the sequences $(q_{n,k})$ defined by

$$q_{n,k} = (1 - n^{-1}) q_{n-1,k} + \frac{n^{-1}}{2} (q_{n-2,k-\lambda_2} + q_{n-2,k-\lambda_3}). \quad (16)$$

After that, we find a relation between the sequences $(q_{n,k})$ and the original sequence $(p_{n,k})$, which terminates the proof.

Remark that if $(q_{n_0,k})_k$ and $(q_{n_0+1,k})_k$ define two probability distributions (i.e. $\sum_k q_{n_0,k} = \sum_k q_{n_0+1,k} = 1$, and $q_{n_0,k}$ and $q_{n_0+1,k}$ are non-negative for every k), then by a simple induction, $(q_{n,k})_k$ also defines a probability distribution for *all* integers $n \geq n_0$. In this case, we can define for every $n \geq n_0$ a random variable Y_n such that $\mathbb{P}(Y_n = k) = q_{n,k}$. The following lemma states that Y_n tends to a Gaussian law.

Lemma 20. *Set $n_0 \geq 0$. Let us consider $(q_{n,k})_{n \geq n_0, k \geq 0}$ a sequence of numbers that:*

- *defines a probability distribution,*

- satisfies (16) after $n \geq n_0 + 2$,
- has a finite support when $n = n_0$ and $n = n_0 + 1$ (that is, the number of k such that $q_{n,k} \neq 0$ is finite).

If Y_n denotes the random variable defined as $\mathbb{P}(Y_n = k) = q_{n,k}$, then $\frac{Y_n - \mu \ln(n)}{\sigma \sqrt{\ln(n)}}$ converges in distribution to a standard Gaussian law, where $\mu = \lambda_2/2 + \lambda_3/2$ and $\sigma^2 = \lambda_2^2/2 + \lambda_3^2/2$.

The subtlety of this lemma lies in the fact that we consider sequences $(q_{n,k})$ that are only defined after some fixed number n_0 , without any initial condition. This flexibility on n_0 will be crucial for the final proof.

Proof. We show here that the generating function of the numbers $q_{n,k}$ satisfies the hypotheses of the Quasi-Powers theorem (Theorem 7).

Step 1: completing the sequence. We first complete the sequence $(q_{n,k})$ so that it satisfies (16) for every $n \geq 0$. To do so, we define $q_{n,k}$ for $n \in \{0, \dots, n_0 - 1\}$ by considering (16) as a backward recurrence:

$$q_{n-2,k-\lambda_3} = 2n q_{n,k} + 2(1-n) q_{n-1,k} - q_{n-2,k-\lambda_2},$$

with initial condition $q_{n,0} = 0$ for every $n \in \{0, \dots, n_0 - 1\}$. (We have assumed that λ_3 is smaller than λ_2 . If it is not the case, we can still swap the roles of λ_2 and λ_3 .) Some subtleties appear here. First, the sequence $(q_{n,k})_{n \geq 0}$ thus completed can now take negative values. This will force us to go back to the non-completed probability sequence $(q_{n,k})_{n \geq n_0}$ to use the Quasi-Powers theorem. Secondly, if $\lambda_2 \neq \lambda_3$, the support is not necessarily finite any more. It will add some difficulty to prove the analyticity in u required by the Quasi-Powers theorem, which justifies the next item.

Step 2: proving the analyticity of the coefficients. We show here that $q_n(u)$, defined as $\sum_{k \geq 0} q_{n,k} u^k$, is analytic at $u = 1$ for every $n \leq n_0 + 1$. It holds for $n = n_0$ and $n = n_0 + 1$ because by assumption, $q_n(u)$ is a polynomial (the support of $(q_{n,k})_k$ is finite). For the numbers smaller than n_0 , we show by induction on $\ell \in \{0, \dots, n_0\}$ that the quantity $r_\ell(u)$, defined as $(u^{\lambda_2} + u^{\lambda_3})^\ell q_{n_0-\ell}(u)$, is a polynomial. Indeed, we observe that (16) can be written as

$$(u^{\lambda_2} + u^{\lambda_3}) q_{n-2}(u) = 2n q_n(u) + 2(1-n) q_{n-1}(u),$$

which implies for every $\ell \geq 0$ by multiplying both sides by $(u^{\lambda_2} + u^{\lambda_3})^{\ell-1}$:

$$r_\ell(u) = 2(n_0 - \ell + 2) (u^{\lambda_2} + u^{\lambda_3}) r_{\ell-2}(u) - 2(n_0 - \ell + 1) r_{\ell-1}(u).$$

The last equality shows that the induction hypothesis is preserved, so by induction (the base case $\ell = 0$ and $\ell = 1$ are obvious), the series $r_\ell(u)$ is polynomial for every $\ell \in \{0, \dots, n_0\}$. In particular, it means that $q_n(u)$ is analytic at 1 for every $0 \leq n \leq n_0$.

Step 3: solving the differential equation. Let us consider $Q(z, u) := \sum_{n,k \geq 0} q_{n,k} z^n u^k$ the completed generating function of the numbers $q_{n,k}$. In the same spirit as the proof of Proposition 12, Equation (16) can be translated in terms of a differential equation on $Q(z, u)$:

$$(1-z) \frac{\partial Q}{\partial z}(z, u) - \frac{z}{2} (u^{\lambda_2} + u^{\lambda_3}) Q(z, u) = q_1(u),$$

where $q_1(u)$ is the coefficient of z^1 in $Q(z, u)$. This equation has for solution

$$Q(z, u) = q_0(u) e^{-\alpha(u)z} (1-z)^{-\alpha(u)} + q_1(u) e^{-\alpha(u)z} (1-z)^{-\alpha(u)} \int_0^z e^{\alpha(u)x} (1-x)^{\alpha(u)-1} dx, \quad (17)$$

where $\alpha(u) := (u^{\lambda_2} + u^{\lambda_3})/2$ and $q_0(u)$ is the constant coefficient in z of $Q(z, u)$.

Step 4: finding a representation of $Q(z, u)$. Let us prove that $Q(z, u)$ has a representation of the form

$$Q(z, u) = A(z, u) + B(z, u) (1 - z)^{-\alpha(u)}.$$

For that, we observe by a simple calculation that an antiderivative for $e^{\alpha(u)z} (1 - z)^{\alpha(u)-1}$ is given by

$$-e^{\alpha(u)} \sum_{k \geq 0} \frac{(1 - z)^{\alpha(u)+k}}{(\alpha(u) + k)k!} (-\alpha(u))^k.$$

Thus, if $H(z, u)$ denotes the series $\sum_{k \geq 0} z^k / ((u + k)k!)$ (which is analytic for every z and for $u = 1$), then (17) can be put into the form

$$\begin{aligned} Q(z, u) = & -q_1(u) e^{\alpha(u)(1-z)} H(\alpha(u)z - \alpha(u), \alpha(u)) \\ & + \left(q_0(u) + q_1(u) e^{\alpha(u)} H(-\alpha(u), \alpha(u)) \right) e^{-\alpha(u)z} (1 - z)^{\alpha(u)}, \end{aligned}$$

which is exactly the wanted representation.

Step 5: application of the Quasi-Powers Theorem. As announced in Step 1, we use the Quasi-Powers Theorem (Theorem 7) not on $Q(z, u)$ (since it could have negative coefficients), but on the non-completed probability generating function $\sum_{n \geq n_0} q_{n,k} z^n u^k$. The latter function has a representation of the form $A(z, u) + B(z, u) (1 - z)^{-\alpha(u)}$ since it differs from $Q(z, u)$ by an analytic function which is here $\sum_{n=0}^{n_0-1} q_n(u) z^n$. Moreover, it is analytic at $(0, 0)$ and it has non-negative coefficients. The variability condition is also satisfied since $\alpha'(1) + \alpha''(1) = \lambda_2^2/2 + \lambda_3^2/2$. The Quasi-Power Theorem thus proves that Y_n converges to a Gaussian limit law with the announced properties. \square

The next lemma, which is quite technical, shows how the error in $O(n^{-2})$ from (15) is propagated over the differences $p_{n,k} - q_{n,k}$, when n goes to $+\infty$.

Lemma 21. *There exists a family of constants $M_{n,k,\ell}$ with $n_0 + 2 \leq \ell \leq n$ and a constant M such that*

- for every $n \geq \ell \geq n_0 + 2$,

$$\sum_{k \geq 0} M_{n,k,\ell} \leq M;$$

- for every $n \geq n_0$ and k positive,

$$|p_{n,k} - q_{n,k}| \leq \max_{k \geq 0} |p_{n_0,k} - q_{n_0,k}| + \max_{k \geq 0} |p_{n_0+1,k} - q_{n_0+1,k}| + \sum_{\ell=n_0+2}^n \frac{M_{n,k,\ell}}{\ell^2};$$

where $p_{n,k}$ is defined by Lemma 19 and $q_{n,k}$ can be any sequence defined by Lemma 20. (In other words, M and the constants $M_{n,k,\ell}$ do not depend on the sequence $q_{n,k}$.)

Proof. The lemma is proved by an induction on n . The statement is obvious for $n = n_0$ and $n = n_0 + 1$.

For $n > n_0 + 1$, the combination of (15) and (16) leads to the inequality

$$\begin{aligned} |p_{n,k} - q_{n,k}| \leq & (1 - n^{-1}) |p_{n-1,k} - q_{n-1,k}| + \frac{n^{-1}}{2} |p_{n-2,k-\lambda_2} - q_{n-2,k-\lambda_2}| \\ & + \frac{n^{-1}}{2} |p_{n-2,k-\lambda_3} - q_{n-2,k-\lambda_3}| + O(n^{-2}). \end{aligned}$$

Referring to the proof of Lemma 19, we see that the error $O(n^{-2})$ in (15) corresponds to the probability $\mathbb{P}\left(X_n \notin \bigcup_{i=1,2,3} \mathcal{C}_i \text{ and } f(X_n) = k\right)$. This number is also $\mu_{n,k} \times \mathbb{P}\left(X_n \notin \bigcup_{i=1,2,3} \mathcal{C}_i\right)$,

where $\mu_{n,k}$ is the conditional probability

$$\mu_{n,k} := \mathbb{P} \left(f(X_n) = k \mid X_n \notin \bigcup_{i=1,2,3} C_i \right).$$

(We have $\sum_{k \geq 0} \mu_{n,k} = 1$.) We have already stated that $\mathbb{P} \left(X_n \notin \bigcup_{i=1,2,3} C_i \right) = O(n^{-2})$, so there exists a constant M such that $\mathbb{P} \left(X_n \notin \bigcup_{i=1,2,3} C_i \right)$ is smaller than $M n^{-2}$ for every $n \geq 0$. Using that fact and the induction hypothesis, the previous inequality becomes

$$\begin{aligned} |p_{n,k} - q_{n,k}| \leq (1 - n^{-1}) \left(x + \sum_{\ell=n_0+2}^{n-1} \frac{M_{n-1,k,\ell}}{\ell^2} \right) + \frac{n^{-1}}{2} \left(x + \sum_{\ell=n_0+2}^{n-2} \frac{M_{n-2,k-\lambda_2,\ell}}{\ell^2} \right) \\ + \frac{n^{-1}}{2} \left(x + \sum_{\ell=n_0+2}^{n-2} \frac{M_{n-2,k-\lambda_3,\ell}}{\ell^2} \right) + M \mu_{n,k} n^{-2}, \end{aligned}$$

where $x := \max_{k \geq 0} |p_{n_0,k} - q_{n_0,k}| + \max_{k \geq 0} |p_{n_0+1,k} - q_{n_0+1,k}|$. Reorganising the terms, we find that

$$|p_{n,k} - q_{n,k}| \leq x + \sum_{\ell=n_0+2}^n \frac{M_{n,k,\ell}}{\ell^2},$$

where we have set

$$M_{n,k,n} = M \mu_{n,k}, \quad M_{n,k,n-1} = (1 - n^{-1}) M_{n-1,k,n-1},$$

and for $n_0 + 2 \leq \ell \leq n - 2$,

$$M_{n,k,\ell} = (1 - n^{-1}) M_{n-1,k,\ell} + \frac{n^{-1}}{2} M_{n-2,k-\lambda_2,\ell} + \frac{n^{-1}}{2} M_{n-2,k-\lambda_3,\ell}.$$

We have $\sum_{k \geq 0} M_{n,k,n} = M$ since $\sum_{k \geq 0} \mu_{n,k} = 1$. As for $n_0 + 2 \leq \ell \leq n - 1$, the induction hypothesis shows that

$$\sum_{k \geq 0} M_{n,k,\ell} \leq (1 - n^{-1}) M + \frac{n^{-1}}{2} M + \frac{n^{-1}}{2} M = M.$$

(The change of variable $k \leftarrow k - \lambda_i$ implies $\sum_{k \geq 0} M_{n-2,k-\lambda_i,\ell} = \sum_{k \geq 0} M_{n-2,k,\ell} \leq M$.) The induction is thus proved. \square

We now have all the tools we need to show Theorem 16.

Proof of Theorem 16. We want to prove that $\frac{f(X_n) - \mu \ln(n)}{\sigma \sqrt{\ln(n)}}$ converges in distribution to a standard Gaussian law, that is, for every $\varepsilon > 0$ and every real number t , there exists $n_1 \geq 0$ such that for every $n \geq n_1$,

$$\left| \mathbb{P} \left(\frac{f(X_n) - \mu \ln(n)}{\sigma \sqrt{\ln(n)}} \leq t \right) - F_{\mathcal{N}}(t) \right| \leq \varepsilon,$$

where $F_{\mathcal{N}}(t)$ denotes the cumulative distribution function of the standard Gaussian law.

1. Definition of n_0 . The series $\sum_{\ell \geq 1} \ell^{-2}$ is convergent, hence its remainder tends to 0. So there exists a number n_0 such that for every $n \geq n_0 + 2$,

$$\sum_{\ell=n_0+2}^n \frac{M}{\ell^2} \leq \sum_{\ell=n_0+2}^{+\infty} \frac{M}{\ell^2} \leq \frac{\varepsilon}{2}, \quad (18)$$

where M is the constant defined by Lemma 21.

2. Definition of an adapted sequence $(q_{n,k})$. Let us define a sequence $(q_{n,k})$ satisfying (16) with initial conditions $q_{n_0,k} := p_{n_0,k}$ and $q_{n_0+1,k} := p_{n_0+1,k}$ for every $k \geq 0$. (The sequence $(q_{n,k})$ satisfies Lemma 20. The finiteness of the support comes from the fact there cannot be more integers k such that $p_{n,k} \neq 0$ than the number of connected diagrams of size n .)

3. Definition of n_1 and first piece of the inequality. We know by Lemma 20 that the variable $\frac{Y_n - \mu \ln(n)}{\sigma \sqrt{\ln(n)}}$ converges in distribution to the standard Gaussian law. So there exists $n_1 \geq n_0 + 2$ such that for every $n \geq n_1$,

$$\left| \mathbb{P} \left(\frac{Y_n - \mu \ln(n)}{\sigma \sqrt{\ln(n)}} \leq t \right) - F_{\mathcal{N}}(t) \right| \leq \frac{\varepsilon}{2}. \quad (19)$$

4. Second piece of the inequality. By definition, we have for every $n \geq n_1$,

$$\begin{aligned} \left| \mathbb{P} \left(\frac{f(X_n) - \mu \ln(n)}{\sigma \sqrt{\ln(n)}} \leq t \right) - \mathbb{P} \left(\frac{Y_n - \mu \ln(n)}{\sigma \sqrt{\ln(n)}} \leq t \right) \right| \\ = \left| \sum_{0 \leq k \leq \mu \ln(n) + t \sigma \sqrt{\ln(n)}} (p_{n,k} - q_{n,k}) \right| \leq \sum_{k \geq 0} |p_{n,k} - q_{n,k}|. \end{aligned}$$

Lemma 21 yields an upper bound for the difference $p_{n,k} - q_{n,k}$, so the previous number is bounded by

$$\sum_{k \geq 0} \left(\max_{k \geq 0} |p_{n_0,k} - q_{n_0,k}| + \max_{k \geq 0} |p_{n_0+1,k} - q_{n_0+1,k}| + \sum_{\ell=n_0+2}^n \frac{M_{n,k,\ell}}{\ell^2} \right).$$

However, by definition of $(q_{n,k})$, $\max_{k \geq 0} |p_{n_0,k} - q_{n_0,k}| = \max_{k \geq 0} |p_{n_0+1,k} - q_{n_0+1,k}| = 0$. We can then swap the sum over k and the sum over ℓ , and use the condition $\sum_{k \geq 0} M_{n,k,\ell} \leq M$ from Lemma (21) to obtain

$$\left| \mathbb{P} \left(\frac{f(X_n) - \mu \ln(n)}{\sigma \sqrt{\ln(n)}} \leq t \right) - \mathbb{P} \left(\frac{Y_n - \mu \ln(n)}{\sigma \sqrt{\ln(n)}} \leq t \right) \right| \leq \sum_{\ell=n_0+2}^n \frac{M}{\ell^2} \leq \frac{\varepsilon}{2}, \quad (20)$$

where the last inequality comes from (18).

5. Conclusion. The conjunction of (19) and (20) shows via a triangle inequality that

$$\left| \mathbb{P} \left(\frac{f(X_n) - \mu \ln(n)}{\sigma \sqrt{\ln(n)}} \leq t \right) - F_{\mathcal{N}}(t) \right| \leq \frac{\varepsilon}{2} + \frac{\varepsilon}{2} = \varepsilon$$

for every $n \geq n_1$, as we had to prove. \square

4.3. POSITION OF THE FIRST TERMINAL CHORD

In this subsection, we are interested by the average position of the first terminal chord for the intersection order. This parameter is relevant since it appears in the sum (2) characterizing the Green function solution of (1).

As an introductory remark, note that the first terminal chord is *always* the chord with the rightmost endpoint, as stated by the following proposition.

Proposition 22. *For every connected diagram, the first terminal chord is the chord that contains the last point of the diagram.*

Proof. We proceed by induction on the number of chords. The property obviously holds when there is only one chord. Assuming now that there are several chords in the diagram, we remove the root chord from the diagram, which creates one or several connected components. By definition of the intersection order, the chords in the topmost component, which we denote C_1 , are smaller than the other ones. But because it is the topmost component, C_1 must also contain the chord with the rightmost endpoint. Therefore, by using the induction hypothesis, the latter chord is the first terminal chord of C_1 , hence the first terminal chord of the whole original diagram. \square

Now let us turn on f_n , the random variable that returns the position of the first terminal chord, under the uniform distribution on connected chord diagrams of size n .

Note that f_n does not satisfy the hypotheses of Theorem 16. Indeed, we can observe that the position of the first terminal chord for every diagram in \mathcal{C}_2 is 2, regardless of the position of the first terminal chord of C_2 .

This remark can be checked experimentally; the observed limit law is not Gaussian. In fact, it seems that f_n/n converges to a law with a density, as shown by Figure 5. We think that this density is $(1 - s)^{-1/2}/2$, with $s \in [0, 1)$. To our knowledge, such a limit law has never been observed on a class of combinatorial objects. This should be the subject of future work.

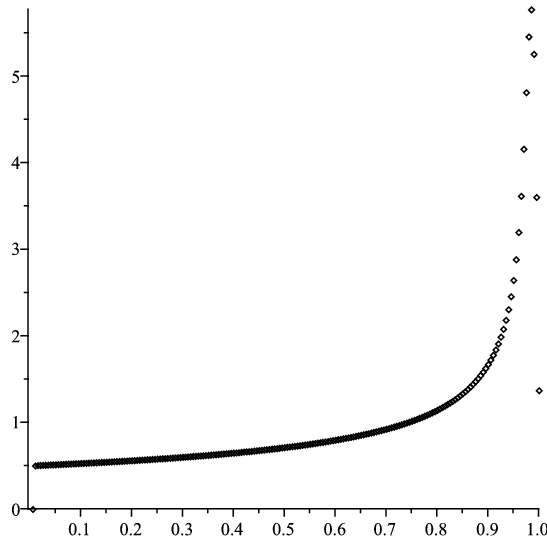


FIGURE 5. Graph of $n \times \mathbb{P}(f_n/n = k)$ in terms of k , for $n = 200$. If f_n have a local limit law as we expect (cf [8, p. 695] for a definition of *local limit law*), then $n \times \mathbb{P}(x \leq f_n/n \leq x + dx)$ should converge to the density of the limit law of f_n/n at the point x .

We calculate here the expected value of this limit law.

Theorem 23. *Let f_n be the position of the first terminal chord of a uniformly distributed random connected chord diagram of size n . The expected value of f_n is asymptotically equivalent to $\frac{2n}{3}$.*

Once again, the proof is based on the approximation of f_n by another law which is easier to study. If we denote by $b(C)$ the position of the first terminal chord of a connected diagram C , and by $\mathcal{C}_1, \mathcal{C}_2, \mathcal{C}_3$ the three sets of connected diagrams defined in the beginning of this section, we can see that

- $b(C) = b(C_1) + 1$ for $C \in \mathcal{C}_1$;
- $b(C) = 2$ for $C \in \mathcal{C}_2$;
- $b(C) = b(C_3) + 1$ for $C \in \mathcal{C}_3$.

A direct adaptation of the proof of Lemma 19 shows then for $n \geq 2$, and $k \in \{3, \dots, n\}$

$$\mathbb{P}(f_n = k) = \left(1 - \frac{1}{n}\right) \mathbb{P}(f_{n-1} = k-1) + \frac{1}{2n} \mathbb{P}(f_{n-2} = k-1) + O(n^{-2}), \quad (21)$$

and

$$\mathbb{P}(f_n = 2) = \frac{1}{2n} + O(n^{-2}).$$

We then define the numbers $g_{n,k}$ thanks to the recurrence

$$g_{n,k} := \left(1 - \frac{1}{n}\right) g_{n-1,k-1} + \frac{1}{2n} g_{n-2,k-1}, \quad g_{n,2} := \frac{1}{2n} \quad (22)$$

for $n \geq 4$, and with initial conditions $g_{n,k} := \mathbb{P}(f_n = k)$ for $n \in \{1, 2, 3\}$. Remark by a straightforward induction that $\sum_{k \geq 0} g_{n,k} = 1$ for every integer n .

We start by proving that the expected values of f_n and $g_{n,k}$ coincide asymptotically.

Lemma 24. *We have*

$$\mathbb{E}(f_n) - \sum_{k \geq 0} k g_{n,k} = o(n).$$

Proof. Set $\varepsilon_n := \mathbb{E}(f_n) - \sum_{k \geq 0} k g_{n,k}$. For $n \geq 2$, the law f_n and the numbers $g_{n,k}$ have $\{2, \dots, n\}$ as a support, hence $\varepsilon_n = \sum_{k=2}^n k (\mathbb{P}(f_n = k) - g_{n,k})$.

Using Equations (21) and (22), we can then deduce that $\varepsilon_n = e_1 + e_2 + e_3$ for $n \geq 4$, where

$$\begin{aligned} e_1 &:= \left(1 - \frac{1}{n}\right) \sum_{k=3}^n k (\mathbb{P}(f_{n-1} = k-1) - g_{n-1,k-1}) \\ e_2 &:= \frac{1}{2n} \sum_{k=3}^n k (\mathbb{P}(f_{n-2} = k-1) - g_{n-2,k-1}), \\ e_3 &:= 2(\mathbb{P}(f_n = 2) - g_{n,2}) + \sum_{k=2}^n O(n^{-2}) = O(n^{-1}), \end{aligned}$$

But we can note that

$$\begin{aligned} e_1 \left(1 - \frac{1}{n}\right)^{-1} &= \sum_{k=3}^n (k-1) (\mathbb{P}(f_{n-1} = k-1) - g_{n-1,k-1}) + \sum_{k=3}^n \mathbb{P}(f_{n-1} = k-1) - \sum_{k=3}^n g_{n-1,k-1} \\ &= \varepsilon_{n-1} + 1 - 1 = \varepsilon_{n-1}. \end{aligned}$$

Similarly, $e_2 = \frac{1}{2n} \varepsilon_{n-2}$ so that

$$\varepsilon_n = \left(1 - \frac{1}{n}\right) \varepsilon_{n-1} + \frac{1}{2n} \varepsilon_{n-2} + O(n^{-1}).$$

The sequence ε_n/n is bounded (because $-n \leq f_n - g_n \leq n$), so has a limit point, let us say, ℓ . We have then

$$\ell = \frac{(n-1)^2}{n^2} \ell + \frac{n-2}{2n^2} \ell + O(n^{-2}).$$

The right-side member is asymptotically equivalent to $\ell - 3\ell/(2n)$. We must then have $\ell = 0$ so that this asymptotic estimate coincides with ℓ . Consequently, the sequence ε_n/n is bounded and have only one limit point, which is 0. Thus ε_n/n tends to 0, which means that $\mathbb{E}(f_n) - \sum_{k \geq 0} k g_{n,k} = o(n)$. \square

The next step is the explicit calculation of the generating function of the numbers $g_{n,k}$.

Lemma 25. *The ordinary generating function $G(z, u)$ of the numbers $g_{n,k}$, namely $\sum_{n,k \geq 0} g_{n,k} z^n u^k$, is equal to*

$$G(z, u) = e^{-\frac{z}{2}} (1 - uz)^{-\frac{1}{2u}} \int_0^z P(x, u) \frac{e^{\frac{x}{2}}}{1-x} (1-ux)^{\frac{1}{2u}-1} dx, \quad (23)$$

where $P(x, u) = (1-x)(u + xu^2 + x^2u^2/4 + x^2u^3/4) + x^3u^2/2$.

Proof. Using (22), we can check that $G(z, u)$ satisfies the linear differential equation

$$uzG(z, u) + 2(1-uz) \frac{\partial G}{\partial z}(z, u) = \frac{u^2 z^2}{1-z} + 2u + 2u^2 z + \frac{1}{2} u^2 z^2 - \frac{1}{2} z^2 u^3.$$

We then verify that this differential equation is solved by (23). \square

Proof of Theorem 23. The sum $\sum_{k \geq 0} k g_{n,k}$ is the n th coefficient of the series $\frac{\partial G}{\partial u}(z, 1)$, where G is the generating function defined in Lemma 25. We are going to use the transfer theorem on $\frac{\partial G}{\partial u}(z, 1)$. This series does not have a non-integral expression, but it is still possible to compute its singular expansion.

Write $G(z, u) = h_1(z, u) \int_0^z h_2(x, u) dx$, where

$$h_1(z, u) = e^{-\frac{z}{2}} (1-uz)^{-\frac{1}{2u}}, \quad h_2(x, u) = P(x, u) \frac{e^{\frac{x}{2}}}{1-x} (1-ux)^{\frac{1}{2u}-1}.$$

We have

$$\frac{\partial G}{\partial u}(z, 1) = \frac{\partial h_1}{\partial u}(z, 1) \int_0^z h_2(x, 1) dx + h_1(z, 1) \int_0^z \frac{\partial h_2}{\partial u}(x, 1) dx.$$

(Since we only have analytic functions, integration and differentiation with respect to u are swappable.) One can explicitly compute the first part:

$$\frac{\partial h_1}{\partial u}(z, 1) \int_0^z h_2(x, 1) dx = \frac{z^2}{2} (1-z)^{-2} + \frac{z}{2} \ln(1-z) (1-z)^{-1},$$

which is asymptotically equivalent to $(1-z)^{-2}/2$ when z approaches 1. Concerning the second part, we calculate $\frac{\partial h_2}{\partial u}(x, 1)$ and observe that

$$\frac{\partial h_2}{\partial u}(x, 1) \sim \frac{e^{\frac{1}{2}}}{4} (1-x)^{-5/2}.$$

We then use Theorem VI.9 from [8, p. 420] to integrate this expansion:

$$\int_0^z \frac{\partial h_2}{\partial u}(x, 1) dx \sim \frac{e^{\frac{1}{2}}}{6} (1-z)^{-3/2},$$

and hence

$$h_1(z, 1) \int_0^z \frac{\partial h_2}{\partial u}(x, 1) dx \sim \frac{1}{6} (1-z)^{-2}.$$

Finally we have $\frac{\partial G}{\partial u}(z, 1) \sim \frac{2}{3} (1-z)^{-2}$, so by the transfer theorem, we have

$$\sum_{k \geq 0} k g_{n,k} \sim \frac{2}{3} n.$$

We conclude thanks to Lemma 24. \square

5. CONCLUSION

In summary, this document establishes numerous exact and asymptotic results on connected chord diagrams. It shows how to compute the next-to- i -leading log expansions, along with their asymptotic regimes. It also shows the Gaussian behaviour of many variables, like the number of terminal chords, and yields their means.

From a combinatorial point of view, this entire study is interesting on its own. It develops new methods to analyse parameters in a context which is not favourable to analytic combinatorics (*a priori*). Moreover, it displays a non-Gaussian limit law, which seems to be new, and maybe deserves a deeper study.

Looking at this from a physical perspective, we observe the dominance of f_0 and f_1 , which respectively denote the residue and the constant term of the Laurent expansion of the regularized Feynman integral of the one loop graph. This is particularly striking for the next-to- i -leading log expansions, whose asymptotic behaviour is governed by f_0 and f_1 (cf (14)). But this dominance can also be noted to a lesser extent to an unrestricted uniform distribution. In fact, by Corollaries 17 and 18, the numbers f_0 and f_1 are on average exponentiated $n - \ln n$ and $\ln n/2$ times in the monomial $f_0^{|C|-k} \prod_{j=2}^k f_{t_j-t_{j-1}}$, which leaves only $\ln n/2$ extra factors for the other f_i (always on average).

To have more information on these extra factors, it would be interesting to study the distribution of the gaps $t_j - t_{j-1}$ other than 1. Conjecturally the number of j such that $t_j - t_{j-1} = \ell$, where ℓ is fixed, asymptotically behaves like a Gaussian law with a mean and variance proportional to $\ln n/n^{\ell-1}$. It should display a double regime: one is discrete – a gap is equal to 1 with a probability $1/2$; the other is continuous – a gap conditioned to be different from 1 should obey to a continuous limit law with mean $2n/\ln n$. The nature of the variance would be also interesting to know.

As for the next steps, the authors intend to generalize their results in the light of [9]. Specifically, the generalization should concern any number of primitives and Dyson-Schwinger equations of various shapes (including the QED shape).

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